Closed trajectories on symmetric convex Hamiltonian energy surfaces

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Abstract

In this article, let $\Sigma \subset \mathbf{R}^{2n}$ be a compact convex Hamiltonian energy surface which is symmetric with respect to the origin. where $n \geq 2$. We prove that there exist at least two geometrically distinct symmetric closed trajectories of the Reeb vector field on Σ .

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1 Introduction and main results

In this article, let Σ be a fixed C^3 compact convex hypersurface in \mathbf{R}^{2n} , i.e., Σ is the boundary of a compact and strictly convex region U in \mathbf{R}^{2n} . We denote the set of all such hypersurfaces by $\mathcal{H}(2n)$. Without loss of generality, we suppose U contains the origin. We denote the set of all compact convex hypersurfaces which are symmetric with respect to the origin by $\mathcal{SH}(2n)$, i.e., $\Sigma = -\Sigma$ for $\Sigma \in \mathcal{SH}(2n)$. We consider closed characteristics (τ, y) on Σ , which are solutions of the following problem

$$\begin{cases}
\dot{y} = JN_{\Sigma}(y), \\
y(\tau) = y(0),
\end{cases}$$
(1.1)

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where $J = \begin{pmatrix} 0 & -I_n \\ I_n & 0 \end{pmatrix}$, I_n is the identity matrix in \mathbf{R}^n , $\tau > 0$ and $N_{\Sigma}(y)$ is the outward normal vector of Σ at y normalized by the condition $N_{\Sigma}(y) \cdot y = 1$. Here $a \cdot b$ denotes the standard inner product of $a, b \in \mathbf{R}^{2n}$. A closed characteristic (τ, y) is prime if τ is the minimal period of y. Two closed characteristics (τ, y) and (σ, z) are geometrically distinct if $y(\mathbf{R}) \neq z(\mathbf{R})$. We denote by $\mathcal{T}(\Sigma)$ the set of geometrically distinct closed characteristics (τ, y) on Σ . A closed characteristic (τ, y) on $\Sigma \in \mathcal{SH}(2n)$ is symmetric if $\{y(\mathbf{R})\} = \{-y(\mathbf{R})\}$, non-symmetric if $\{y(\mathbf{R})\} \cap \{-y(\mathbf{R})\} = \emptyset$. It was proved in [LLZ] that a prime characteristic (τ, y) on $\Sigma \in \mathcal{SH}(2n)$ is symmetric if and only if $y(t) = -y(t + \frac{\tau}{2})$ for all $t \in \mathbf{R}$.

There is a long standing conjecture on the number of closed characteristics on compact convex hypersurfaces in \mathbb{R}^{2n} :

$$^{\#}\mathcal{T}(\Sigma) \ge n, \qquad \forall \ \Sigma \in \mathcal{H}(2n).$$
 (1.2)

Since the pioneering works [Rab1] of P. Rabinowitz and [Wei1] of A. Weinstein in 1978 on the existence of at least one closed characteristic on every hypersurface in $\mathcal{H}(2n)$, the existence of multiple closed characteristics on $\Sigma \in \mathcal{H}(2n)$ has been deeply studied by many mathematicians. When $n \geq 2$, besides many results under pinching conditions, in 1987-1988 I. Ekeland-L. Lassoued, I. Ekeland-H. Hofer, and A, Szulkin (cf. [EkL1], [EkH1], [Szu1]) proved

$$^{\#}\mathcal{T}(\Sigma) \geq 2, \qquad \forall \, \Sigma \in \mathcal{H}(2n).$$

In [HWZ] of 1998, H. Hofer-K. Wysocki-E. Zehnder proved that ${}^{\#}\mathcal{T}(\Sigma) = 2$ or ∞ holds for every $\Sigma \in \mathcal{H}(4)$. In [LoZ1] of 2002, Y. Long and C. Zhu proved

$$^{\#}\mathcal{T}(\Sigma) \ge \left[\frac{n}{2}\right] + 1, \qquad \forall \, \Sigma \in \mathcal{H}(2n),$$

where we denote by $[a] \equiv \max\{k \in \mathbf{Z} \mid k \leq a\}$. In [WHL], the authors proved the conjecture for n = 3. In [LLZ], the the authors proved the conjecture for $\Sigma \in \mathcal{SH}(2n)$.

Note that in [W2], the author proved if ${}^{\#}\mathcal{T}(\Sigma) = n$ for some $\Sigma \in \mathcal{SH}(2n)$ and n = 2 or 3, then any $(\tau, y) \in \mathcal{T}(\Sigma)$ is symmetric. Thus it is natural to conjecture that

$$^{\#}\mathcal{T}_s(\Sigma) \ge n, \qquad \forall \ \Sigma \in \mathcal{SH}(2n),$$
 (1.3)

where $T_s(\Sigma)$ denotes the set of geometrically distinct symmetric closed characteristics (τ, y) on Σ . The following is the main result in this article:

Theorem 1.1. We have $\#\mathcal{T}_s(\Sigma) \geq 2$ for any $\Sigma \in \mathcal{SH}(2n)$, where $n \geq 2$.

In this article, let N, N_0 , Z, Q, R, and C denote the sets of natural integers, non-negative integers, rational numbers, real numbers, and complex numbers respectively. Denote by

 $a \cdot b$ and |a| the standard inner product and norm in \mathbf{R}^{2n} . Denote by $\langle \cdot, \cdot \rangle$ and $\| \cdot \|$ the standard L^2 -inner product and L^2 -norm. For an S^1 -space X, we denote by X_{S^1} the homotopy quotient of X module the S^1 -action, i.e., $X_{S^1} = S^{\infty} \times_{S^1} X$. We define the functions

$$\begin{cases}
[a] = \max\{k \in \mathbf{Z} \mid k \le a\}, & E(a) = \min\{k \in \mathbf{Z} \mid k \ge a\}, \\
\varphi(a) = E(a) - [a],
\end{cases} (1.4)$$

Specially, $\varphi(a) = 0$ if $a \in \mathbf{Z}$, and $\varphi(a) = 1$ if $a \notin \mathbf{Z}$. In this article we use only **Q**-coefficients for all homological modules.

2 A variational structure for closed characteristics

In this section, we transform the problem (1.1) into a fixed period problem of a Hamiltonian system and then study its variational structure.

In the rest of this paper, we fix a $\Sigma \in \mathcal{SH}(2n)$ and assume the following condition on Σ :

(F) There exist only finitely many geometrically distinct symmetric closed characteristics $\{(\tau_j, y_j)\}_{1 \le j \le k}$ on Σ .

Note that $(\tau, y) \in \mathcal{T}_s(\Sigma)$ is a solution of (1.1) if and only if it satisfies the equation

$$\begin{cases} \dot{y} = JN_{\Sigma}(y), \\ y(\frac{\tau}{2}) = -y(0), \end{cases}$$
 (2.1)

Now we construct a variational structure of closed characteristics as the following.

lemma 2.1. (cf. Proposition 2.2 of [WHL]) For any sufficiently small $\vartheta \in (0,1)$, there exists a function $\varphi \equiv \varphi_{\vartheta} \in C^{\infty}(\mathbf{R}, \mathbf{R}^{+})$ depending on ϑ which has 0 as its unique critical point in $[0, +\infty)$ such that the following hold

- (i) $\varphi(0) = 0 = \varphi'(0)$ and $\varphi''(0) = 1 = \lim_{t \to 0^+} \frac{\varphi'(t)}{t}$.
- (ii) $\varphi(t)$ is a polynomial of degree 2 in a neighborhood of $+\infty$.
- $(iii) \ \tfrac{d}{dt} \left(\tfrac{\varphi'(t)}{t} \right) < 0 \ for \ t > 0, \ and \ \lim_{t \to +\infty} \tfrac{\varphi'(t)}{t} < \vartheta, \ i.e., \ \tfrac{\varphi'(t)}{t} \ is \ strictly \ decreasing \ for \ t > 0.$
- (iv) $\min(\frac{\varphi'(t)}{t}, \varphi''(t)) \ge \sigma$ for all $t \in \mathbf{R}^+$ and some $\sigma > 0$. Consequently, φ is strictly convex on $[0, +\infty)$.
- (v) In particular, we can choose $\alpha \in (1,2)$ sufficiently close to 2 and $c \in (0,1)$ such that $\varphi(t) = ct^{\alpha}$ whenever $\frac{\varphi'(t)}{t} \in [\vartheta, 1 \vartheta]$ and t > 0.
- Let $j: \mathbf{R}^{2n} \to \mathbf{R}$ be the gauge function of Σ , i.e., $j(\lambda x) = \lambda$ for $x \in \Sigma$ and $\lambda \geq 0$, then $j \in C^3(\mathbf{R}^{2n} \setminus \{0\}, \mathbf{R}) \cap C^0(\mathbf{R}^{2n}, \mathbf{R})$ and $\Sigma = j^{-1}(1)$. Denote by $\hat{\tau} = \inf_{1 \leq j \leq k} \tau_j$ and $\hat{\sigma} = \min\{|y|^2 \mid y \in \Sigma\}$.

By the same proof of Proposition 2.4 of [WHL], we have the following

Proposition 2.2. Let $a > \hat{\tau}$, $\vartheta_a \in \left(0, \frac{1}{a} \min\{\hat{\tau}, \hat{\sigma}\}\right)$ and φ_a be a C^{∞} function associated to ϑ_a satisfying (i)-(iv) of Lemma 2.1. Define the Hamiltonian function $H_a(x) = a\varphi_a(j(x))$ and consider the fixed period problem

$$\begin{cases} \dot{x}(t) = JH'_a(x(t)) \\ x(\frac{1}{2}) = -x(0) \end{cases}$$
 (2.2)

Then the following hold:

(i) $H_a \in C^3(\mathbf{R}^{2n} \setminus \{0\}, \mathbf{R}) \cap C^1(\mathbf{R}^{2n}, \mathbf{R})$ and there exist R, r > 0 such that

$$r|\xi|^2 \le H_a''(x)\xi \cdot \xi \le R|\xi|^2$$
, $\forall x \in \mathbf{R}^{2n} \setminus \{0\}, \ \xi \in \mathbf{R}^{2n}$.

(ii) There exist $\epsilon_1, \epsilon_2 \in \left(0, \frac{1}{2}\right)$ and $C \in \mathbf{R}$, such that

$$\frac{\epsilon_1|x|^2}{2} - C \le H_a(x) \le \frac{\epsilon_2|x|^2}{2} + C, \quad \forall x \in \mathbf{R}^{2n}.$$

- (iii) Solutions of (2.2) are $x \equiv 0$ and $x = \rho y(\tau t)$ with $\frac{\varphi_a'(\rho)}{\rho} = \frac{\tau}{a}$, where (τ, y) is a solution of (2.1). In particular, nonzero solutions of (2.2) are in one to one correspondence with solutions of (2.1) with period $\tau < a$.
 - (iv) There exists $r_0 > 0$ independent of a and there exists $\mu_a > 0$ depending on a such that

$$H_a''(x)\xi \cdot \xi \ge 2ar_0|\xi|^2$$
, for $0 < |x| \le \mu_a$, $\xi \in \mathbf{R}^{2n}$.

In the following, we will use the Clarke-Ekeland dual action principle. As usual, the Fenchel transform of a function $F: \mathbf{R}^{2n} \to \mathbf{R}$ is defined by

$$F^*(y) = \sup\{x \cdot y - F(x) \mid x \in \mathbf{R}^{2n}\}. \tag{2.3}$$

Following Proposition 2.2.10 of [Eke3], Lemma 3.1 of [Eke1] and the fact that $F_1 \leq F_2 \Leftrightarrow F_1^* \geq F_2^*$, we have:

Proposition 2.3. Let H_a be a function defined in Proposition 2.2 and $G_a = H_a^*$ the Fenchel transform of H_a . Then we have

(i)
$$G_a \in C^2(\mathbf{R}^{2n} \setminus \{0\}, \mathbf{R}) \cap C^1(\mathbf{R}^{2n}, \mathbf{R})$$
 and

$$G_a'(y) = x \Leftrightarrow y = H_a'(x) \Rightarrow H_a''(x)G_a''(y) = 1.$$

(ii) G_a is strictly convex. Let R and r be the real numbers given by (i) of Proposition 2.2. Then we have

$$R^{-1}|\xi|^2 \le G_a''(y)\xi \cdot \xi \le r^{-1}|\xi|^2, \quad \forall y \in \mathbf{R}^{2n} \setminus \{0\}, \ \xi \in \mathbf{R}^{2n}.$$

(iii) Let $\epsilon_1, \epsilon_2, C$ be the real numbers given by (ii) of Proposition 2.2. Then we have

$$\frac{|x|^2}{2\epsilon_2} - C \le G_a(x) \le \frac{|x|^2}{2\epsilon_1} + C, \quad \forall x \in \mathbf{R}^{2n}.$$

(iv) Let $r_0 > 0$ be the constant given by (iv) of Proposition 2.2. Then there exists $\eta_a > 0$ depending on a such that the following holds

$$G''_a(y)\xi \cdot \xi \le \frac{1}{2ar_0}|\xi|^2$$
, for $0 < |y| \le \eta_a, \ \xi \in \mathbf{R}^{2n}$.

(v) In particular, let $H_a = a\varphi_a(j(x))$ with φ_a satisfying further (v) of Lemma 2.1. Then we have $G_a(\mu j'(z)) = c_1 \mu^{\beta}$ when $z \in \Sigma$ and $\mu j'(z) \in \{H'_a(x) \mid H_a(x) = acj(x)^{\alpha}\}$, where c is given by (v) of Lemma 2.1, $c_1 > 0$ is some constant depending on a and $\alpha^{-1} + \beta^{-1} = 1$.

Now we apply the dual action principle to problem (2.3). Let

$$L^{2}\left(\mathbf{R} / \left(\frac{1}{2}\mathbf{Z}\right), \mathbf{R}^{2n}\right) = \{u \in L^{2}(\mathbf{R}, \mathbf{R}^{2n}) | u(t+1/2) = -u(t)\}.$$
 (2.4)

Define a linear operator $M: L^2\left(\mathbf{R} \left/ \left(\frac{1}{2}\mathbf{Z}\right), \mathbf{R}^{2n}\right) \to L^2\left(\mathbf{R} \left/ \left(\frac{1}{2}\mathbf{Z}\right), \mathbf{R}^{2n}\right)\right.$ by

$$\frac{d}{dt}Mu(t) = u(t). (2.5)$$

Lemma 2.4. M is a compact operator from $L^2\left(\mathbf{R} \left/ \left(\frac{1}{2}\mathbf{Z}\right), \mathbf{R}^{2n}\right)\right)$ into itself and $M^* = -M$.

Proof. Note that M sends $L^2\left(\mathbf{R} / \left(\frac{1}{2}\mathbf{Z}\right), \mathbf{R}^{2n}\right)$ into $W^{1,2}\left([0, 1/2], \mathbf{R}^{2n}\right)$, and the identity map from $W^{1,2}\left([0, 1/2], \mathbf{R}^{2n}\right)$ to $L^2\left(\mathbf{R} / \left(\frac{1}{2}\mathbf{Z}\right), \mathbf{R}^{2n}\right)$ is compact by the Rellich-Kondrachov theorem. Hence M is compact.

To check it is anti-symmetric, we use integrate by parts:

$$\int_0^{1/2} (Mu, v)dt = -\int_0^{1/2} (u, Mv)dt + (Mu, Mv)|_0^{1/2}.$$

and the last term vanishes since Mu(1/2) = -Mu(0) and Mv(1/2) = -Mv(0). Hence M is anti-symmetric.

The dual action functional on $L^{2}\left(\mathbf{R}\left/\left(\frac{1}{2}\mathbf{Z}\right),\mathbf{R}^{2n}\right)$ is defined by

$$\Psi_a(u) = \int_0^{1/2} \left(\frac{1}{2} J u \cdot M u + G_a(-J u) \right) dt, \tag{2.6}$$

where G_a is given by Proposition 2.3.

Proposition 2.5. The functional Ψ_a is bounded from below on $L^2\left(\mathbf{R} / \left(\frac{1}{2}\mathbf{Z}\right), \mathbf{R}^{2n}\right)$.

Proof. For any $u \in L^2\left(\mathbf{R} / \left(\frac{1}{2}\mathbf{Z}\right), \mathbf{R}^{2n}\right)$, we represent u by its Fourier series

$$u(t) = \sum_{k \in 2\mathbf{Z}+1} e^{2k\pi Jt} x_k, \quad x_k \in \mathbf{R}^{2n}.$$

$$(2.7)$$

Then we have

$$Mu(t) = -J \sum_{k \in 2\mathbb{Z}+1} \frac{1}{2\pi k} e^{2k\pi J t} x_k.$$
 (2.8)

Hence

$$\frac{1}{2}\langle Ju, Mu \rangle = -\frac{1}{2} \sum_{k \in 2\mathbb{Z}+1} \frac{1}{2\pi k} |x_k|^2 \ge -\frac{1}{4\pi} ||u||^2.$$
 (2.9)

By (2.6), we have

$$\Psi_{a}(u) = \int_{0}^{1/2} \left(\frac{1}{2}Ju \cdot Mu + G_{a}(-Ju)\right) dt
\geq \frac{1}{2}\langle Ju, Mu \rangle + \int_{0}^{1/2} \left(\frac{|u|^{2}}{2\epsilon_{2}} - C\right) dt.
\geq \left(\frac{1}{2\epsilon_{2}} - \frac{1}{4\pi}\right) ||u||^{2} - C
\geq C_{1}||u||^{2} - C$$
(2.10)

for some constant $C_1 > 0$, where in the first inequality, we have used (iii) of Proposition 2.3. Hence the proposition holds.

Proposition 2.6. The functional Ψ_a is $C^{1,1}$ on $L^2\left(\mathbf{R} / \left(\frac{1}{2}\mathbf{Z}\right), \mathbf{R}^{2n}\right)$ and satisfies the Palais-Smale condition. Suppose x is a solution of (2.2), then $u = \dot{x}$ is a critical point of Ψ_a . Conversely, suppose u is a critical point of Ψ_a , then Mu is a solution of (2.2). In particular, solutions of (2.2) are in one to one correspondence with critical points of Ψ_a .

Proof. By (ii) of Proposition 2.3 and the same proof of Proposition 3.3 on p.33 of [Eke1], we have Ψ_a is $C^{1,1}$ on $L^2\left(\mathbf{R} / \left(\frac{1}{2}\mathbf{Z}\right), \mathbf{R}^{2n}\right)$. By (2.10) and the proof of Lemma 5.2.8 of [Eke3], we have Ψ_a satisfies the Palais-Smale condition.

By (2.6), we have

$$\langle \Psi_a'(u), v \rangle = \langle Mu, Jv \rangle - \langle G_a'(-Ju), Jv \rangle, \tag{2.11}$$

where we use the fact that

$$Mu(t) = \int_0^t u(s)ds - \frac{1}{2} \int_0^{1/2} u(s)ds,$$
 (2.12)

and MJu(t) = JMu(t). Hence $\Psi'_a(u) = 0$ if and only if $Mu = G'_a(-Ju)$, where we used the fact $G'_a(-Ju(\frac{1}{2})) = G'_a(Ju(0)) = -G'_a(-Ju(0))$. Taking Frenchel dual we have $-Ju = H'_a(Mu)$, i.e., $u = JH'_a(Mu)$. Hence Mu is a solution of (2.2). The converse is obvious.

Proposition 2.7. We have $\Psi_a(u_a) < 0$ for every critical point $u_a \neq 0$ of Ψ_a .

Proof. By Propositions 2.2 and 2.6, we have $u_a = \dot{x}_a$ and $x_a = \rho_a y(\tau t)$ with

$$\frac{\varphi_a'(\rho_a)}{\rho_a} = \frac{\tau}{a}. (2.13)$$

Hence we have

$$\Psi_{a}(u_{a}) = \int_{0}^{1/2} \left(\frac{1}{2}J\dot{x}_{a} \cdot x_{a} + G_{a}(-J\dot{x}_{a})\right) dt$$

$$= -\frac{1}{4}\langle H'_{a}(x_{a}), x_{a}\rangle + \int_{0}^{1/2} G_{a}(H'_{a}(x_{a})) dt$$

$$= \frac{1}{4}a\varphi'_{a}(\rho_{a})\rho_{a} - \frac{1}{2}a\varphi_{a}(\rho_{a}).$$
(2.14)

Here the second equality follows from (2.2) and the third equality follows from (i) of Proposition 2.3 and (2.3).

Let $f(t) = \frac{1}{2}a\varphi_a'(t)t - a\varphi_a(t)$ for $t \ge 0$. Then we have f(0) = 0 and $f'(t) = \frac{a}{2}(\varphi_a''(t)t - \varphi_a'(t)) < 0$ since $\frac{d}{dt}(\frac{\varphi_a'(t)}{t}) < 0$ by (iii) of Lemma 2.1. This together with (2.13) yield the proposition.

We have a natural S^1 -action on $L^2\left(\mathbf{R} \left/ \left(\frac{1}{2}\mathbf{Z}\right), \mathbf{R}^{2n}\right)\right)$ defined by

$$\theta * u(t) = u(\theta + t), \quad \forall \theta \in S^1 \equiv \mathbf{R}/\mathbf{Z}, t \in \mathbf{R}.$$
 (2.15)

Then we have

Lemma 2.8. The functional Ψ_a is S^1 -invariant.

Proof. Note that we have the following

Claim. We have $M(\theta * u) = \theta * (Mu)$.

In fact, by (2.12), we have

$$M(\theta * u)(t) = \int_0^t \theta * u(s)ds - \frac{1}{2} \int_0^{1/2} \theta * u(s)ds$$
$$= \int_0^t u(\theta + s)ds - \frac{1}{2} \int_0^{1/2} u(\theta + s)ds$$
$$= \int_\theta^{t+\theta} u(s)ds - \frac{1}{2} \int_\theta^{1/2+\theta} \cdot u(s)ds$$

On the other hand, we have

$$\theta * (Mu)(t) = \theta * \left(\int_0^t u(s)ds - \frac{1}{2} \int_0^{1/2} u(s)ds \right)$$

$$= \int_0^{t+\theta} u(s)ds - \frac{1}{2} \int_0^{1/2} u(s)ds$$

$$= \int_0^{\theta} u(s)ds + \int_{\theta}^{t+\theta} u(s)ds - \frac{1}{2} \int_0^{\theta} u(s)ds - \frac{1}{2} \int_{\theta}^{1/2} u(s)ds$$

$$= \frac{1}{2} \int_0^{\theta} u(s)ds + \int_{\theta}^{t+\theta} u(s)ds - \frac{1}{2} \int_{\theta}^{1/2} u(s)ds$$

$$= -\frac{1}{2} \int_{1/2}^{\theta+1/2} u(s)ds + \int_{\theta}^{t+\theta} u(s)ds - \frac{1}{2} \int_{\theta}^{1/2} u(s)ds$$

$$= \int_{\theta}^{t+\theta} u(s)ds - \frac{1}{2} \int_{\theta}^{1/2+\theta} u(s)ds,$$
(2.16)

where in (2.16), we use the fact u(t+1/2) = -u(t). Hence the claim holds.

Now we have

$$\begin{split} \Psi_{a}(\theta * u) &= \int_{0}^{1/2} \left(\frac{1}{2} J(\theta * u) \cdot M(\theta * u) + G_{a}(-J(\theta * u)) \right) dt, \\ &= \int_{0}^{1/2} \left(\frac{1}{2} \theta * (Ju) \cdot \theta * (Mu) + G_{a}(\theta * (-Ju)) \right) dt, \\ &= \int_{\theta}^{\theta + 1/2} \left(\frac{1}{2} Ju \cdot Mu + G_{a}(-Ju) \right) dt \\ &= \int_{\theta}^{1/2} \left(\frac{1}{2} Ju \cdot Mu + G_{a}(-Ju) \right) dt + \int_{1/2}^{\theta + 1/2} \left(\frac{1}{2} Ju \cdot Mu + G_{a}(-Ju) \right) dt \\ &= \int_{\theta}^{1/2} \left(\frac{1}{2} Ju \cdot Mu + G_{a}(-Ju) \right) dt + \int_{0}^{\theta} \left(\frac{1}{2} (-Ju) \cdot (-Mu) + G_{a}(Ju) \right) dt \\ &= \int_{\theta}^{1/2} \left(\frac{1}{2} Ju \cdot Mu + G_{a}(-Ju) \right) dt = \Psi_{a}(u), \end{split}$$

where in the above computation, we use u(t+1/2) = -u(t) and $G_a(x) = G_a(-x)$, which follows from $\Sigma = -\Sigma$. Hence the proposition holds.

For any $\kappa \in \mathbf{R}$, we denote by

$$\Lambda_a^{\kappa} = \left\{ u \in L^2 \left(\mathbf{R} / \left(\frac{1}{2} \mathbf{Z} \right), \mathbf{R}^{2n} \right) \mid \Psi_a(u) \le \kappa \right\}. \tag{2.17}$$

For a critical point u of Ψ_a , we denote by

$$\Lambda_a(u) = \Lambda_a^{\Psi_a(u)} = \left\{ w \in L^2\left(\mathbf{R} \middle/ \left(\frac{1}{2}\mathbf{Z}\right), \mathbf{R}^{2n}\right) \middle| \Psi_a(w) \le \Psi_a(u) \right\}. \tag{2.18}$$

Clearly, both sets are S^1 -invariant. Since the S^1 -action preserves Ψ_a , if u is a critical point of Ψ_a , then the whole orbit $S^1 \cdot u$ is formed by critical points of Ψ_a . Denote by $crit(\Psi_a)$ the set of critical points of Ψ_a . Note that by the condition (F), (iii) of Proposition 2.2 and Proposition 2.6, the number of critical orbits of Ψ_a is finite. Hence as usual we can make the following definition.

Definition 2.9. Suppose u is a nonzero critical point of Ψ_a , and \mathcal{N} is an S^1 -invariant open neighborhood of $S^1 \cdot u$ such that $crit(\Psi_a) \cap (\Lambda_a(u) \cap \mathcal{N}) = S^1 \cdot u$. Then the S^1 -critical modules of $S^1 \cdot u$ is defined by

$$C_{S^{1}, q}(\Psi_{a}, S^{1} \cdot u) = H_{S^{1}, q}(\Lambda_{a}(u) \cap \mathcal{N}, (\Lambda_{a}(u) \setminus S^{1} \cdot u) \cap \mathcal{N})$$

$$\equiv H_{q}((\Lambda_{a}(u) \cap \mathcal{N})_{S^{1}}, ((\Lambda_{a}(u) \setminus S^{1} \cdot u) \cap \mathcal{N})_{S^{1}}), \tag{2.19}$$

where $H_{S^1,*}$ is the S^1 -equivariant homology with rational coefficients in the sense of A. Borel (cf. Chapter IV of [Bor1]).

By the same argument as Proposition 3.2 of [WHL], we have the following proposition for critical modules.

Proposition 2.10. The critical module $C_{S^1, q}(\Psi_a, S^1 \cdot u)$ is independent of the choice of H_a defined in Proposition 2.2 in the sense that if x_i are solutions of (2.2) with Hamiltonian functions $H_{a_i}(x) \equiv a_i \varphi_{a_i}(j(x))$ for i = 1 and 2 respectively such that both x_1 and x_2 correspond to the same closed characteristic (τ, y) on Σ . Then we have

$$C_{S^1, q}(\Psi_{a_1}, S^1 \cdot \dot{x}_1) \cong C_{S^1, q}(\Psi_{a_2}, S^1 \cdot \dot{x}_2), \quad \forall q \in \mathbf{Z}.$$
 (2.20)

In other words, the critical modules are invariant for all $a > \tau$ and φ_a satisfying (i)-(iv) of Lemma 2.1.

In order to compute the critical modules, as in p.35 of [Eke1] and p.219 of [Eke3] we introduce the following.

Definition 2.11. Suppose u is a nonzero critical point of Ψ_a . Then the formal Hessian of Ψ_a at u is defined by

$$Q_a(v, v) = \int_0^{1/2} (Jv \cdot Mv + G_a''(-Ju)Jv \cdot Jv)dt,$$
 (2.21)

which defines an orthogonal splitting $L^2\left(\mathbf{R} \left/ \left(\frac{1}{2}\mathbf{Z}\right), \mathbf{R}^{2n}\right) = E_- \oplus E_0 \oplus E_+ \text{ of } L^2\left(\mathbf{R} \left/ \left(\frac{1}{2}\mathbf{Z}\right), \mathbf{R}^{2n}\right)\right)$ into negative, zero and positive subspaces. The index of u is defined by $i(u) = \dim E_-$ and the nullity of u is defined by $\nu(u) = \dim E_0$.

Next we show that the index and nullity defined as above are the Morse index and nullity of a corresponding functional on a finite dimensional subspace of $L^2\left(\mathbf{R} \left/ \left(\frac{1}{2}\mathbf{Z}\right), \mathbf{R}^{2n}\right)\right)$.

- **Lemma 2.12.** Let Ψ_a be a functionals defined by (2.6). Then there exists a finite dimensional S^1 -invariant subspace X of $L^2\left(\mathbf{R} / \left(\frac{1}{2}\mathbf{Z}\right), \mathbf{R}^{2n}\right)$ and a S^1 -equivariant map $h_a: X \to X^{\perp}$ such that the following hold
 - (i) For $g \in X$, each function $h \mapsto \Psi_a(g+h)$ has $h_a(g)$ as the unique minimum in X^{\perp} . Let $\psi_a(g) = \Psi_a(g+h_a(g))$. Then we have
- (ii) The function ψ_a is C^1 on X and S^1 -invariant. g_a is a critical point of ψ_a if and only if $g_a + h_a(g_a)$ is a critical point of Ψ_a .
- (iii) If $g_a \in X$ and H_a is C^k with $k \geq 2$ in a neighborhood of the trajectory of $g_a + h_a(g_a)$, then ψ_a is C^{k-1} in a neighborhood of g_a . In particular, if g_a is a nonzero critical point of ψ_a , then ψ_a is C^2 in a neighborhood of the critical orbit $S^1 \cdot g_a$. The index and nullity of Ψ_a at $g_a + h_a(g_a)$ defined in Definition 2.11 coincide with the Morse index and nullity of ψ_a at g_a .

(iv) For any $\kappa \in \mathbf{R}$, we denote by

$$\widetilde{\Lambda}_a^{\kappa} = \{ g \in X \mid \psi_a(g) \le \kappa \}. \tag{2.22}$$

Then the natural embedding $\widetilde{\Lambda}_a^{\kappa} \hookrightarrow \Lambda_a^{\kappa}$ given by $g \mapsto g + h_a(g)$ is an S^1 -equivariant homotopy equivalence.

Proof. By (ii) of Proposition 2.3, we have

$$(G_a'(u) - G_a'(v), u - v) \ge \omega |u - v|^2, \quad \forall u, v \in \mathbf{R}^{2n},$$
 (2.23)

for some $\omega > 0$. Hence we can use the proof of Proposition 3.9 of [Vit1] to obtain X and h_a . In fact, X is the subspace of $L^2\left(\mathbf{R} / \left(\frac{1}{2}\mathbf{Z}\right), \mathbf{R}^{2n}\right)$ generated by the eigenvectors of -JM whose eigenvalues are less than $-\frac{\omega}{2}$ and $h_a(g)$ is defined by the equation

$$\frac{\partial}{\partial h}\Psi_a(g + h_a(g)) = 0,$$

then (i)-(iii) follows from Proposition 3.9 of [Vit1]. (iv) follows from Lemma 5.1 of [Vit1].

Note that Ψ_a is not C^2 in general, and then we can not apply Morse theory to Ψ_a directly. After the finite dimensional approximation, the function ψ_a has much better differentiability, which allows us to apply the Morse theory to study its property.

Proposition 2.13. Let Ψ_a be a functional defined by (2.6), and $u_a = \dot{x}_a$ be the critical point of Ψ_a so that x_a corresponds to a closed characteristic (τ, y) on Σ . Then the nullity $\nu(u_a)$ of the functional Ψ_a at its critical point u_a is the number of linearly independent solutions of the boundary value problem

$$\begin{cases} \dot{\xi}(t) = JH_a''(x_a(t))\xi \\ \xi(\frac{1}{2}) = -\xi(0) \end{cases}$$
 (2.24)

Proof. By (2.21), we have

$$Q_a(v, w) = \int_0^{1/2} (Jv \cdot Mw + G_a''(-Ju)Jv \cdot Jw)dt,$$

= $\langle Mw, Jv \rangle + \langle (H_a''(x_a(t))^{-1}Jw, Jv \rangle$ (2.25)

where we have used (2.2) and (i) of Proposition 2.3. Now $w \in E_0$ if and only if $Q_a(v, w) = 0$ for any $v \in L^2\left(\mathbf{R} / \left(\frac{1}{2}\mathbf{Z}\right), \mathbf{R}^{2n}\right)$. Hence we must have $Mw + (H''_a(x_a(t))^{-1}Jw = 0$, i.e., we have $w = JH''_a(x_a(t))Mw$. Hence Mw solves (2.25).

Denote by R(t) the fundamental solution of the linearized system

$$\dot{\xi}(t) = JH_a''(x_a(t))\xi(t), \tag{2.26}$$

Then we have the following

Proposition 2.14. In an appropriate coordinates there holds

$$R(1/2) = \begin{pmatrix} A & B \\ 0 & C \end{pmatrix}$$
 with $A = \begin{pmatrix} -1 & -\gamma \\ 0 & -1 \end{pmatrix}$,

with $\gamma > 0$ and C is independent of H_a .

Proof. Note that by Lemma 1.6.11 of [Eke3], we have

$$R(t)T_{y(0)}\Sigma \subset T_{y(\tau t)}\Sigma.$$
 (2.27)

Differentiating (2.2) and use the fact $x_a(t+1/2) = -x_a(t)$, we have

$$R(1/2)\dot{x}_a(0) = -\dot{x}_a(0). \tag{2.28}$$

Let

$$x_a(\rho, t) = \rho y \left(\frac{\tau t}{T_\rho}\right) \quad \text{with } \frac{\tau}{T_\rho} = \frac{a\varphi_a'(\rho)}{\rho}.$$
 (2.29)

Then we have $x_a(\rho, T_{\rho}/2) = -x_a(\rho, 0)$. Differentiating it with respect to ρ and using (2.29) together with $\dot{x}_a(1/2) = -\dot{x}_a(0)$, we get

$$-\frac{\tau}{2a}\frac{d}{d\rho}\left(\frac{\rho}{\varphi_a'(\rho)}\right)\dot{x}_a(0) + R(1/2)\rho^{-1}x_a(0) = -\rho^{-1}x_a(0).$$

Hence we have

$$R(1/2)x_a(0) = -x_a(0) + \frac{\rho\tau}{2a} \frac{d}{d\rho} \left(\frac{\rho}{\varphi_a'(\rho)}\right) \dot{x}_a(0) = x_a(0) + \gamma \dot{x}_a(0), \tag{2.30}$$

where $\gamma > 0$ since $\frac{d}{d\rho} \left(\frac{\rho}{\varphi_a'(\rho)} \right) > 0$ by (iii) of Proposition 2.1. For any $w \in \mathbf{R}^{2n}$, we have

$$H_a''(x_a)w = a\varphi_a''(j(x_a))(j'(x_a), w)j'(x_a) + a\varphi_a'(j(x_a))j''(x_a)w$$

$$= a\varphi_a''(j(x_a))(j'(y), w)j'(y) + \tau j''(y)w.$$
(2.31)

The last equality follows from (iii) of Proposition 2.2. Let z(t) = R(t)z(0) for $z(0) \in T_{y(0)}\Sigma$. Then by (2.27), we have $\dot{z}(t) = \tau j''(y(t))z(t)$. Therefore $R(1/2)|_{T_{y(0)}}\Sigma$ is independent of the choice of H_a in Proposition 2.2. Summing up, we have proved that in an appropriate coordinates there holds

$$R(1/2) = \begin{pmatrix} A & B \\ 0 & C \end{pmatrix}$$
 with $A = \begin{pmatrix} -1 & -\gamma \\ 0 & -1 \end{pmatrix}$,

with C is independent of H_a , where we use $\{-\dot{x}_a(0), x_a(0), e_1, \dots, e_{2n-2}\}$ as an basis of \mathbf{R}^{2n} .

Proposition 2.15. Let Ψ_a be a functional defined by (2.6), and u be a nonzero critical point of Ψ_a . Then we have

$$C_{S^1, q}(\Psi_a, S^1 \cdot u) = 0, \qquad \forall q \notin [i(u), i(u) + \nu(u) - 1].$$
 (2.32)

Proof. By (iv) of Lemma 2.12, we have

$$C_{S^1, q}(\Psi_a, S^1 \cdot u) \simeq C_{S^1, q}(\psi_a, S^1 \cdot u),$$
 (2.33)

where $C_{S^1, q}(\psi_a, S^1 \cdot u) = H_{S^1, q}(\widetilde{\Lambda}_a(u) \cap \mathcal{N}, (\widetilde{\Lambda}_a(u) \setminus S^1 \cdot u) \cap \mathcal{N})$ and \mathcal{N} is an S^1 -invariant open neighborhood of $S^1 \cdot u$ such that $crit(\psi_a) \cap (\widetilde{\Lambda}_a(u) \cap \mathcal{N}) = S^1 \cdot u$. By (iii) of Lemma 2.12, the functional ψ_a is C^2 near $S^1 \cdot u$. Thus we can use the Gromoll-Meyer theory in the equivariant sense to obtain the proposition.

Recall that for a principal U(1)-bundle $E \to B$, the Fadell-Rabinowitz index (cf. [FaR1]) of E is defined to be $\sup\{k \mid c_1(E)^{k-1} \neq 0\}$, where $c_1(E) \in H^2(B, \mathbf{Q})$ is the first rational Chern class. For a U(1)-space, i.e., a topological space X with a U(1)-action, the Fadell-Rabinowitz index is defined to be the index of the bundle $X \times S^{\infty} \to X \times_{U(1)} S^{\infty}$, where $S^{\infty} \to CP^{\infty}$ is the universal U(1)-bundle. For any $\kappa \in \mathbf{R}$, we denote by

$$\Psi_a^{\kappa-} = \left\{ w \in L^2 \left(\mathbf{R} / \left(\frac{1}{2} \mathbf{Z} \right), \mathbf{R}^{2n} \right) \mid \Psi_a(w) < \kappa \right\}. \tag{2.34}$$

Then as in P.218 of [Eke3], we define

$$c_i = \inf\{\delta \in \mathbf{R} \mid \hat{I}(\Psi_a^{\kappa-}) \ge i\},\tag{2.35}$$

where \hat{I} is the Fadell-Rabinowitz index given above. Then as Proposition 3 in P.218 of [Eke3], we have

Proposition 2.16. Every c_i is a critical value of Ψ_a . If $c_i = c_j$ for some i < j, then there are infinitely many geometrically distinct symmetric closed characteristics on Σ .

By a similar argument as Proposition 3.5 of [W1] and Proposition 2.15, we have

Proposition 2.17. Suppose u is the critical point of Ψ_a found in Proposition 2.16. Then we have

$$\Psi_a(u) = c_i, \qquad C_{S^1, 2(i-1)}(\Psi_a, S^1 \cdot u) \neq 0.$$
 (2.36)

In particular, we have $i(u) \le 2(i-1) \le i(u) + \nu(u) - 1$.

3 Index iteration theory for symmetric closed characteristics

In this section, we study the index iteration theory for symmetric closed characteristics.

Note that if $(\tau, y) \in \mathcal{T}_s(\Sigma)$, then $((2m-1)\tau, y)$ is a solution of (2.1) for any $m \in \mathbb{N}$. Thus $((2m-1)\tau, y)$ corresponds to a critical point of Ψ_a via Propositions 2.2 and 2.6, we denote it by u^{2m-1} . First note that we have the following

Lemma 3.1. Suppose u^{2m-1} is a nonzero critical point of Ψ_a such that u corresponds to $(\tau, y) \in \mathcal{T}_s(\Sigma)$. Let $H(x) = j(x)^2$, where j is the gauge function of Σ . Then $i(u^{2m-1})$ equals the index of the following quadratic form

$$Q_{(2m-1)\tau/2}(\xi, \xi) = \int_0^{(2m-1)\tau/2} (J\dot{\xi} \cdot \xi + (H''(y(t)))^{-1} J\dot{\xi} \cdot J\dot{\xi}) dt, \tag{3.1}$$

where $\xi \in W^{1,2}\left(\mathbf{R} / \left(\frac{(2m-1)\tau}{2}\mathbf{Z}\right), \mathbf{R}^{2n}\right) \equiv \{w \in W^{1,2}(\mathbf{R}, \mathbf{R}^{2n}) | w(t + \frac{(2m-1)\tau}{2}) = -w(t)\}.$ Moreover, we have $\nu(u^{2m-1}) = \text{nullity}Q_{(2m-1)\tau/2} - 1$.

Proof. By a similar argument as in proposition 1.7.5 and P.36 of [Eke3] and Proposition 3.5 of [WHL], we obtain the lemma.

Suppose u^{2k-1} is a nonzero critical point of Ψ_a such that u corresponds to $(\tau, y) \in \mathcal{T}_s(\Sigma)$. Then for any $\omega \in \mathbf{U}$, let

$$Q_{(2k-1)\tau/2}^{\omega}(\xi, \xi) = \int_{0}^{(2k-1)\tau/2} (J\dot{\xi} \cdot \xi + (H''(y(t)))^{-1} J\dot{\xi} \cdot J\dot{\xi}) dt, \tag{3.2}$$

where $\xi \in E^{\omega}_{(2k-1)\tau/2} \equiv \{u \in W^{1,2}([0,(2k-1)\tau/2],\mathbf{C}^{2n})|w(\frac{(2k-1)\tau}{2}) = \omega w(0)\}..$

Clearly the quadratic form $Q_{(2m-1)\tau/2}$ on the real Hilbert space $W^{1,2}\left(\mathbf{R} \left/ \left(\frac{(2m-1)\tau}{2}\mathbf{Z}\right),\mathbf{R}^{2n}\right)\right)$ and the Hermitian form $Q_{(2m-1)\tau/2}^{-1}$ on the complex Hilbert space $E_{(2m-1)\tau/2}^{-1}$ have the same index. If $\omega^{2m-1}=-1$, we identify $E_{\tau/2}^{\omega}$ with a subspace of $E_{(2m-1)\tau/2}^{-1}$ via

$$E_{\tau/2}^{\omega} = \{ u \in W^{1,2}(\mathbf{R}, \mathbf{C}^{2n}) | w(t + \tau/2) = \omega w(t) \}.$$
(3.3)

Note that if $\xi \in E^{\omega}_{\tau/2}$, we have

$$Q_{(2m-1)\tau/2}^{\omega}(\xi, \, \xi) = \int_{0}^{(2m-1)\tau/2} (J\dot{\xi} \cdot \xi + (H''(y(t)))^{-1} J\dot{\xi} \cdot J\dot{\xi}) dt$$

$$= \sum_{k=0}^{2m-1} (\omega \overline{\omega})^{k} \int_{0}^{\tau/2} (J\dot{\xi} \cdot \xi + (H''(y(t)))^{-1} J\dot{\xi} \cdot J\dot{\xi}) dt$$

$$= (2m-1)Q_{\tau/2}^{\omega}(\xi, \xi). \tag{3.4}$$

Lemma 3.2. The spaces $E_{\tau/2}^{\omega}$ for $\omega^{2m-1} = -1$ are orthogonal subspaces of $E_{(2m-1)\tau/2}^{-1}$, both for the standard Hilbert structure and for $Q_{(2m-1)\tau/2}^{-1}$, and we have the decomposition

$$E_{(2m-1)\tau/2}^{-1} = \bigoplus_{\omega^{2m-1} = -1} E_{\tau/2}^{\omega}.$$
 (3.5)

Proof. Any $\xi \in E^{-1}_{(2m-1)\tau/2}$ can be written as

$$\xi(t) = \sum_{p \in 2\mathbf{Z}+1} x_p \exp\left(\frac{2i\pi pt}{(2m-1)\tau}\right)$$
(3.6)

for q = 1, 3, ..., 4m - 3, denote by C(q) the set of all p such that $p - q \in (4m - 2)\mathbf{Z}$. Thus we may write

$$\xi(t) = \sum_{\substack{q \in 2\mathbf{Z}+1\\1 \le q \le 4m-3}} \xi_q(t), \qquad \xi_q(t) = \sum_{C(q)} x_p \exp\left(\frac{2i\pi pt}{(2m-1)\tau}\right)$$
(3.7)

Then we have

$$\xi_q(t+\tau/2) = \sum_{C(q)} x_p \exp\left(\frac{2i\pi pt}{(2m-1)\tau} + \frac{i\pi p}{2m-1}\right)$$
$$= \exp\left(\frac{i\pi q}{2m-1}\right) \xi_q(t). \tag{3.8}$$

Thus $\xi_q \in E_{\tau/2}^{\omega}$ with $\omega = \exp\left(\frac{i\pi q}{2m-1}\right)$, when q runs from $1, 3, \ldots, 4m-3$, then ω runs through the 2m-1 roots of -1.

For $\xi \in E^{\omega}_{\tau/2}$ and $\eta \in E^{\lambda}_{\tau/2}$ with $\omega \neq \lambda$ are 2m-1 roots of -1, we have

$$Q_{(2m-1)\tau/2}^{-1}(\xi, \eta) = \int_{0}^{(2m-1)\tau/2} (J\dot{\xi} \cdot \eta + (H''(y(t)))^{-1} J\dot{\xi} \cdot J\dot{\eta}) dt$$

$$= \sum_{k=0}^{2m-1} (\omega \overline{\lambda})^{k} \int_{0}^{\tau/2} (J\dot{\xi} \cdot \eta + (H''(y(t)))^{-1} J\dot{\xi} \cdot J\dot{\eta}) dt$$

$$= 0. \tag{3.9}$$

Thus the lemma holds.

Definition 3.3. We define the Bott maps $j_{\tau/2}$ and $n_{\tau/a}$ from U to Z by

$$j_{\tau/2}(\omega) = \text{index}Q_{\tau/2}^{\omega}, \qquad n_{\tau/2}(\omega) = \text{nullity}Q_{\tau/2}^{\omega},$$
 (3.10)

By Lemmas 3.1 and 3.2, we have

Proposition 3.4. Suppose u^{2m-1} is a nonzero critical point of Ψ_a such that u corresponds to $(\tau, y) \in \mathcal{T}_s(\Sigma)$. Then we have

$$i(u^{2m-1}) = \sum_{\omega^{2m-1} = -1} j_{\tau/2}(\omega) \qquad \nu(u^{2m-1}) = \sum_{\omega^{2m-1} = -1} n_{\tau/2}(\omega) - 1. \tag{3.11}$$

Note that $j_{\tau/2}(\omega)$ coincide with the function defined in Definition 1.5.3 of [Eke3] for the linear Hamiltonian system

$$\begin{cases} \dot{\xi}(t) = JA(t)\xi\\ A(t+\tau/2) = A(t) \end{cases}$$
(3.12)

where A(t) = H''(y(t)). Denote by $i^E(A, k)$ and $\nu^E(A, k)$ the index and nullity of the k-th iteration of the system (3.12) defined by Ekeland in [Eke3]. Denote by i(A, k) and $\nu(A, k)$ the Maslov-type

index and nullity of the k-th iteration of the system (3.12) defined by Conley, Zehnder and Long (cf. §5.4 of [Lon4]). Then we have

Theorem 3.5. (cf. Theorem 15.1.1 of [Lon4]) We have

$$i^{E}(A,k) = i(A,k) - n, \quad \nu^{E}(A,k) = \nu(A,k),$$
 (3.13)

for any $k \in \mathbf{N}$.

Theorem 3.6. Suppose u^{2m-1} is a nonzero critical point of Ψ_a such that u corresponds to $(\tau, y) \in \mathcal{T}_s(\Sigma)$. Then we have

$$i(u^{2m-1}) = i_{-1}(A, 2m-1), \qquad \nu(u^{2m-1}) = \nu_{-1}(A, 2m-1) - 1.$$
 (3.14)

where $i_{-1}(A, k)$ and $\nu_{-1}(A, k)$ are the Maslov-type index and nullity introduced in [Lon2].

Proof. By Corollary 1.5.4 of [Eke3] and Theorem 9.2.1 of [Lon4] respectively, we have

$$i^{E}(A, 4m-2) = i^{E}(A, 2m-1) + i^{E}_{-1}(A, 2m-1),$$

 $i(A, 4m-2) = i(A, 2m-1) + i_{-1}(A, 2m-1)$ (3.15)

and by Lemma 3.1, we have $i(u^{2m-1}) = i_{-1}^{E}(A, 2m-1)$. Thus the theorem follows from Theorem 3.5.

Now we compute $i(u^{2m-1})$ via the index iteration method in [Lon4]. First we recall briefly an index theory for symplectic paths. All the details can be found in [Lon4].

As usual, the symplectic group Sp(2n) is defined by

$$\operatorname{Sp}(2n) = \{ M \in \operatorname{GL}(2n, \mathbf{R}) \, | \, M^T J M = J \},$$

whose topology is induced from that of \mathbf{R}^{4n^2} . For $\tau > 0$ we are interested in paths in $\mathrm{Sp}(2n)$:

$$\mathcal{P}_{\tau}(2n) = \{ \gamma \in C([0,\tau], \operatorname{Sp}(2n)) \, | \, \gamma(0) = I_{2n} \},$$

which is equipped with the topology induced from that of Sp(2n). The following real function was introduced in [Lon2]:

$$D_{\omega}(M) = (-1)^{n-1} \overline{\omega}^n \det(M - \omega I_{2n}), \quad \forall \omega \in \mathbf{U}, M \in \operatorname{Sp}(2n).$$

Thus for any $\omega \in \mathbf{U}$ the following codimension 1 hypersurface in $\mathrm{Sp}(2n)$ is defined in [Lon2]:

$$\operatorname{Sp}(2n)^{0}_{\omega} = \{ M \in \operatorname{Sp}(2n) \mid D_{\omega}(M) = 0 \}.$$

For any $M \in \operatorname{Sp}(2n)^0_{\omega}$, we define a co-orientation of $\operatorname{Sp}(2n)^0_{\omega}$ at M by the positive direction $\frac{d}{dt}Me^{t\epsilon J}|_{t=0}$ of the path $Me^{t\epsilon J}$ with $0 \le t \le 1$ and $\epsilon > 0$ being sufficiently small. Let

$$\operatorname{Sp}(2n)_{\omega}^{*} = \operatorname{Sp}(2n) \setminus \operatorname{Sp}(2n)_{\omega}^{0},$$

$$\mathcal{P}_{\tau,\omega}^{*}(2n) = \{ \gamma \in \mathcal{P}_{\tau}(2n) \mid \gamma(\tau) \in \operatorname{Sp}(2n)_{\omega}^{*} \},$$

$$\mathcal{P}_{\tau,\omega}^{0}(2n) = \mathcal{P}_{\tau}(2n) \setminus \mathcal{P}_{\tau,\omega}^{*}(2n).$$

For any two continuous arcs ξ and $\eta:[0,\tau]\to \operatorname{Sp}(2n)$ with $\xi(\tau)=\eta(0)$, it is defined as usual:

$$\eta * \xi(t) = \begin{cases} \xi(2t), & \text{if } 0 \le t \le \tau/2, \\ \eta(2t - \tau), & \text{if } \tau/2 \le t \le \tau. \end{cases}$$

Given any two $2m_k \times 2m_k$ matrices of square block form $M_k = \begin{pmatrix} A_k & B_k \\ C_k & D_k \end{pmatrix}$ with k = 1, 2, as in [Lon4], the \diamond -product of M_1 and M_2 is defined by the following $2(m_1 + m_2) \times 2(m_1 + m_2)$ matrix $M_1 \diamond M_2$:

$$M_1 \diamond M_2 = \begin{pmatrix} A_1 & 0 & B_1 & 0 \\ 0 & A_2 & 0 & B_2 \\ C_1 & 0 & D_1 & 0 \\ 0 & C_2 & 0 & D_2 \end{pmatrix}.$$

Denote by $M^{\diamond k}$ the k-fold \diamond -product $M\diamond\cdots\diamond M$. Note that the \diamond -product of any two symplectic matrices is symplectic. For any two paths $\gamma_j\in\mathcal{P}_{\tau}(2n_j)$ with j=0 and 1, let $\gamma_0\diamond\gamma_1(t)=\gamma_0(t)\diamond\gamma_1(t)$ for all $t\in[0,\tau]$.

A special path $\xi_n \in \mathcal{P}_{\tau}(2n)$ is defined by

$$\xi_n(t) = \begin{pmatrix} 2 - \frac{t}{\tau} & 0\\ 0 & (2 - \frac{t}{\tau})^{-1} \end{pmatrix}^{\diamond n} \quad \text{for } 0 \le t \le \tau.$$
 (3.16)

Definition 3.7. (cf. [Lon2], [Lon4]) For any $\omega \in \mathbf{U}$ and $M \in \mathrm{Sp}(2n)$, define

$$\nu_{\omega}(M) = \dim_{\mathbf{C}} \ker_{\mathbf{C}}(M - \omega I_{2n}). \tag{3.17}$$

For any $\tau > 0$ and $\gamma \in \mathcal{P}_{\tau}(2n)$, define

$$\nu_{\omega}(\gamma) = \nu_{\omega}(\gamma(\tau)). \tag{3.18}$$

If $\gamma \in \mathcal{P}_{\tau,\omega}^*(2n)$, define

$$i_{\omega}(\gamma) = [\operatorname{Sp}(2n)^{0}_{\omega} : \gamma * \xi_{n}], \tag{3.19}$$

where the right hand side of (3.19) is the usual homotopy intersection number, and the orientation of $\gamma * \xi_n$ is its positive time direction under homotopy with fixed end points.

If $\gamma \in \mathcal{P}^0_{\tau,\omega}(2n)$, we let $\mathcal{F}(\gamma)$ be the set of all open neighborhoods of γ in $\mathcal{P}_{\tau}(2n)$, and define

$$i_{\omega}(\gamma) = \sup_{U \in \mathcal{F}(\gamma)} \inf\{i_{\omega}(\beta) \mid \beta \in U \cap \mathcal{P}_{\tau,\omega}^{*}(2n)\}.$$
 (3.20)

Then

$$(i_{\omega}(\gamma), \nu_{\omega}(\gamma)) \in \mathbf{Z} \times \{0, 1, \dots, 2n\},\$$

is called the index function of γ at ω .

For any $M \in \operatorname{Sp}(2n)$ and $\omega \in \mathbf{U}$, the splitting numbers $S_M^{\pm}(\omega)$ of M at ω are defined by

$$S_M^{\pm}(\omega) = \lim_{\epsilon \to 0^+} i_{\omega \exp(\pm \sqrt{-1}\epsilon)}(\gamma) - i_{\omega}(\gamma), \tag{3.21}$$

for any path $\gamma \in \mathcal{P}_{\tau}(2n)$ satisfying $\gamma(\tau) = M$.

Let $\Omega^0(M)$ be the path connected component containing $M = \gamma(\tau)$ of the set

$$\Omega(M) = \{ N \in \operatorname{Sp}(2n) \mid \sigma(N) \cap \mathbf{U} = \sigma(M) \cap \mathbf{U} \text{ and}$$

$$\nu_{\lambda}(N) = \nu_{\lambda}(M) \ \forall \ \lambda \in \sigma(M) \cap \mathbf{U} \}.$$
(3.22)

Here $\Omega^0(M)$ is called the homotopy component of M in Sp(2n).

In [Lon2]-[Lon4], the following symplectic matrices were introduced as basic normal forms:

$$D(\lambda) = \begin{pmatrix} \lambda & 0 \\ 0 & \lambda^{-1} \end{pmatrix}, \qquad \lambda = \pm 2, \tag{3.23}$$

$$N_1(\lambda, b) = \begin{pmatrix} \lambda & b \\ 0 & \lambda \end{pmatrix}, \qquad \lambda = \pm 1, b = \pm 1, 0, \tag{3.24}$$

$$N_{1}(\lambda, b) = \begin{pmatrix} \lambda & b \\ 0 & \lambda \end{pmatrix}, \qquad \lambda = \pm 1, b = \pm 1, 0,$$

$$R(\theta) = \begin{pmatrix} \cos \theta & -\sin \theta \\ \sin \theta & \cos \theta \end{pmatrix}, \qquad \theta \in (0, \pi) \cup (\pi, 2\pi),$$

$$N_{2}(\omega, b) = \begin{pmatrix} R(\theta) & b \\ 0 & R(\theta) \end{pmatrix}, \qquad \theta \in (0, \pi) \cup (\pi, 2\pi),$$

$$(3.24)$$

$$N_2(\omega, b) = \begin{pmatrix} R(\theta) & b \\ 0 & R(\theta) \end{pmatrix}, \qquad \theta \in (0, \pi) \cup (\pi, 2\pi), \tag{3.26}$$

where $b = \begin{pmatrix} b_1 & b_2 \\ b_3 & b_4 \end{pmatrix}$ with $b_i \in \mathbf{R}$ and $b_2 \neq b_3$.

Splitting numbers possess the following properties:

Lemma 3.8. (cf. [Lon2] and Lemma 9.1.5 of [Lon4]) Splitting numbers $S_M^{\pm}(\omega)$ are well defined, i.e., they are independent of the choice of the path $\gamma \in \mathcal{P}_{\tau}(2n)$ satisfying $\gamma(\tau) = M$ appeared in (3.21). For $\omega \in \mathbf{U}$ and $M \in \mathrm{Sp}(2n)$, splitting numbers $S_N^{\pm}(\omega)$ are constant for all $N \in \Omega^0(M)$. Moreover, we have

$$S_M^{\pm}(\omega) = 0, \quad if \quad \omega \notin \sigma(M).$$

$$S_M^+(\omega) = S_M^-(\overline{\omega}), \quad \forall \omega \in \mathbf{U}.$$

Lemma 3.9. (cf. [Lon2], Lemma 9.1.5 of [Lon4]) For any $M_i \in \operatorname{Sp}(2n_i)$ with i = 0 and 1, there holds

$$S_{M_0 \diamond M_1}^{\pm}(\omega) = S_{M_0}^{\pm}(\omega) + S_{M_1}^{\pm}(\omega), \qquad \forall \ \omega \in \mathbf{U}.$$

$$(3.27)$$

We have the following

Theorem 3.10. (cf. [Lon3] and Theorem 1.8.10 of [Lon4]) For any $M \in \operatorname{Sp}(2n)$, there is a path $f:[0,1] \to \Omega^0(M)$ such that f(0) = M and

$$f(1) = M_1 \diamond \dots \diamond M_l, \tag{3.28}$$

where each M_i is a basic normal form listed in (3.23)-(3.26) for $1 \le i \le l$.

Now we deduce the index iteration formula for each case in (3.23)-(3.26), Note that the splitting numbers are computed in List 9.1.12 of [Lon4].

Case 1. M is conjugate to a matrix $\begin{pmatrix} 1 & b \\ 0 & 1 \end{pmatrix}$ for some b > 0.

In this case, we have $(S_M^+(1), S_M^-(1)) = (1, 1)$. Thus by Theorem 9.2.1 of [Lon4], we have

$$i_{-1}(\gamma^{2m-1}) = \sum_{\omega^{2m-1}=-1} i_{\omega}(\gamma) = \sum_{k=1}^{2m-1} i_{\frac{(2k-1)\pi}{2m-1}}(\gamma) = (2m-1)(i_1(\gamma)+1),$$

$$\nu_{-1}(\gamma^{2m-1}) = 0.$$
(3.29)

Case 2. $M = I_2$, the 2×2 identity matrix.

In this case, we have $(S_M^+(1), S_M^-(1)) = (1, 1)$. Thus as in Case 1, we have

$$i_{-1}(\gamma^{2m-1}) = (2m-1)(i_1(\gamma)+1), \quad \nu_{-1}(\gamma^{2m-1}) = 0.$$
 (3.30)

Case 3. M is conjugate to a matrix $\begin{pmatrix} 1 & b \\ 0 & 1 \end{pmatrix}$ for some b < 0.

In this case, we have $(S_M^+(1), S_M^-(1)) = (0, 0)$. Thus by Theorem 9.2.1 of [Lon4], we have

$$i_{-1}(\gamma^{2m-1}) = \sum_{\omega^{2m-1}=-1} i_{\omega}(\gamma) = \sum_{k=1}^{2m-1} i_{\frac{(2k-1)\pi}{2m-1}}(\gamma) = (2m-1)i_{1}(\gamma),$$

$$\nu_{-1}(\gamma^{2m-1}) = 0.$$
(3.31)

Case 4. M is conjugate to a matrix $\begin{pmatrix} -1 & b \\ 0 & -1 \end{pmatrix}$ for some b < 0.

In this case, we have $(S_M^+(-1), S_M^-(-1)) = (1, 1)$. Thus by Theorem 9.2.1 of [Lon4], we have

$$i_{-1}(\gamma^{2m-1}) = \sum_{\omega^{2m-1} = -1} i_{\omega}(\gamma) = \sum_{k=1}^{2m-1} i_{\frac{(2k-1)\pi}{2m-1}}(\gamma)$$

$$= \sum_{k=1}^{m-1} i_{\frac{(2k-1)\pi}{2m-1}}(\gamma) + i_{-1}(\gamma) + \sum_{k=m+1}^{2m-1} i_{\frac{(2k-1)\pi}{2m-1}}(\gamma)$$

$$= (m-1)i_{1}(\gamma) + i_{1}(\gamma) - 1 + (m-1)(i_{1}(\gamma) - 1 + 1)$$

$$= (2m-1)i_{1}(\gamma) - 1,$$

$$\nu_{-1}(\gamma^{2m-1}) = 1.$$
(3.32)

Case 5. $M = -I_2$.

In this case, we have $(S_M^+(-1), S_M^-(-1)) = (1, 1)$. Thus as in Case 4, we have

$$i_{-1}(\gamma^{2m-1}) = (2m-1)i_1(\gamma) - 1, \quad \nu_{-1}(\gamma^{2m-1}) = 2.$$
 (3.33)

Case 6. M is conjugate to a matrix $\begin{pmatrix} -1 & b \\ 0 & -1 \end{pmatrix}$ for some b > 0.

In this case, we have $(S_M^+(-1), S_M^-(-1)) = (0,0)$. Thus by Theorem 9.2.1 of [Lon4], we have

$$i_{-1}(\gamma^{2m-1}) = \sum_{\omega^{2m-1}=-1} i_{\omega}(\gamma) = \sum_{k=1}^{2m-1} i_{\frac{(2k-1)\pi}{2m-1}}(\gamma) = (2m-1)i_{1}(\gamma),$$

$$\nu_{-1}(\gamma^{2m-1}) = 1.$$
(3.34)

Case 7. $M = \begin{pmatrix} \cos \theta & -\sin \theta \\ \sin \theta & \cos \theta \end{pmatrix}$ with some $\theta \in (0, \pi) \cup (\pi, 2\pi)$.

In this case, we have $(S_M^+(e^{\sqrt{-1}\theta}), S_M^-(e^{\sqrt{-1}\theta})) = (0,1)$. Thus by Theorem 9.2.1 of [Lon4] and Lemma 3.8, we have

$$i_{-1}(\gamma^{2m-1}) = \sum_{\omega^{2m-1}=-1} i_{\omega}(\gamma) = \sum_{k=1}^{2m-1} i_{\frac{(2k-1)\pi}{2m-1}}(\gamma)$$

$$= \sum_{2k-1 < \frac{(2m-1)\theta}{\pi}} i_{1}(\gamma) + \sum_{\frac{(2m-1)\theta}{\pi} \le 2k-1 \le \frac{(2m-1)(2\pi-\theta)}{\pi}} (i_{1}(\gamma) - 1)$$

$$+ \sum_{\frac{(2m-1)(2\pi-\theta)}{\pi} < 2k-1 \le 4m-2} i_{1}(\gamma)$$

$$= (2m-1)(i_{1}(\gamma) - 1) + 2E\left(\frac{(2m-1)\theta}{2\pi} + \frac{1}{2}\right) - 2,$$

$$\nu_{-1}(\gamma^{2m-1}) = 2 - 2\phi\left(\frac{(2m-1)\theta}{2\pi} + \frac{1}{2}\right), \tag{3.35}$$

provided $\theta \in (0, \pi)$. When $\theta \in (\pi, 2\pi)$, we have

$$i_{-1}(\gamma^{2m-1}) = \sum_{\omega^{2m-1}=-1} i_{\omega}(\gamma) = \sum_{k=1}^{2m-1} i_{\frac{(2k-1)\pi}{2m-1}}(\gamma)$$

$$= \sum_{2k-1 \le \frac{(2m-1)(2\pi-\theta)}{\pi}} i_{1}(\gamma) + \sum_{\frac{(2m-1)(2\pi-\theta)}{\pi} < 2k-1 < \frac{(2m-1)\theta}{\pi}} (i_{1}(\gamma)+1)$$

$$+ \sum_{\frac{(2m-1)\theta}{\pi} \le 2k-1 \le 4m-2} i_1(\gamma)$$

$$= (2m-1)(i_1(\gamma)-1) + 2E\left(\frac{(2m-1)\theta}{2\pi} + \frac{1}{2}\right) - 2e^{-1}(\gamma^{2m-1})$$

$$= 2 - 2\phi\left(\frac{(2m-1)\theta}{2\pi} + \frac{1}{2}\right).$$

Case 8. $M = \begin{pmatrix} R(\theta) & b \\ 0 & R(\theta) \end{pmatrix}$ with some $\theta \in (0, \pi) \cup (\pi, 2\pi)$ and $b = \begin{pmatrix} b_1 & b_2 \\ b_3 & b_4 \end{pmatrix} \in \mathbf{R}^{2 \times 2}$, such that $(b_2 - b_3) \sin \theta < 0$.

In this case, we have $(S_M^+(e^{\sqrt{-1}\theta}), S_M^-(e^{\sqrt{-1}\theta})) = (1,1)$. Thus by Theorem 9.2.1 of [Lon4], we have

$$i_{-1}(\gamma^{2m-1}) = \sum_{\omega^{2m-1}=-1} i_{\omega}(\gamma) = \sum_{k=1}^{2m-1} i_{\frac{(2k-1)\pi}{2m-1}}(\gamma)$$

$$= (2m-1)i_{1}(\gamma) + 2\phi \left(\frac{(2m-1)\theta}{2\pi} + \frac{1}{2}\right) - 2,$$

$$\nu_{-1}(\gamma^{2m-1}) = 2 - 2\phi \left(\frac{(2m-1)\theta}{2\pi} + \frac{1}{2}\right). \tag{3.36}$$

Case 9. $M = \begin{pmatrix} R(\theta) & b \\ 0 & R(\theta) \end{pmatrix}$ with some $\theta \in (0, \pi) \cup (\pi, 2\pi)$ and $b = \begin{pmatrix} b_1 & b_2 \\ b_3 & b_4 \end{pmatrix} \in \mathbf{R}^{2 \times 2}$, such hat $(b_2 - b_3) \sin \theta > 0$.

In this case, we have $(S_M^+(e^{\sqrt{-1}\theta}), S_M^-(e^{\sqrt{-1}\theta})) = (0,0)$. Thus by Theorem 9.2.1 of [Lon4], we have

$$i_{-1}(\gamma^{2m-1}) = \sum_{\omega^{2m-1}=-1} i_{\omega}(\gamma) = \sum_{k=1}^{2m-1} i_{\frac{(2k-1)\pi}{2m-1}}(\gamma) = (2m-1)i_{1}(\gamma),$$

$$\nu_{-1}(\gamma^{2m-1}) = 2 - 2\phi \left(\frac{(2m-1)\theta}{2\pi} + \frac{1}{2}\right). \tag{3.37}$$

Case 10. M is hyperbolic, i.e., $\sigma(M) \cap \mathbf{U} = \emptyset$.

In this case, by Theorem 9.2.1 of [Lon4], we have

$$i_{-1}(\gamma^{2m-1}) = \sum_{\omega^{2m-1}=-1} i_{\omega}(\gamma) = \sum_{k=1}^{2m-1} i_{\frac{(2k-1)\pi}{2m-1}}(\gamma) = (2m-1)i_1(\gamma),$$

$$\nu_{-1}(\gamma^{2m-1}) = 0.$$
(3.38)

Proposition 3.11. For any $m \in \mathbb{N}$, we have the estimate

$$i_{-1}(\gamma^{2m+1}) - i_{-1}(\gamma^{2m-1}) - \nu_{-1}(\gamma^{2m-1}) \ge 2i_1(\gamma) - e(M).$$
 (3.39)

Proof. We consider each of the above cases.

Case 1. M is conjugate to a matrix $\begin{pmatrix} 1 & b \\ 0 & 1 \end{pmatrix}$ for some b > 0. In this case we have

$$i_{-1}(\gamma^{2m+1}) - i_{-1}(\gamma^{2m-1}) - \nu_{-1}(\gamma^{2m-1}) = 2i_1(\gamma) + 2.$$

Case 2. $M = I_2$, the 2×2 identity matrix.

In this case we have

$$i_{-1}(\gamma^{2m+1}) - i_{-1}(\gamma^{2m-1}) - \nu_{-1}(\gamma^{2m-1}) = 2i_1(\gamma) + 2.$$

Case 3. M is conjugate to a matrix $\begin{pmatrix} 1 & b \\ 0 & 1 \end{pmatrix}$ for some b < 0.

In this case we have

$$i_{-1}(\gamma^{2m+1}) - i_{-1}(\gamma^{2m-1}) - \nu_{-1}(\gamma^{2m-1}) = 2i_1(\gamma).$$

Case 4. M is conjugate to a matrix $\begin{pmatrix} -1 & b \\ 0 & -1 \end{pmatrix}$ for some b < 0.

In this case we have

$$i_{-1}(\gamma^{2m+1}) - i_{-1}(\gamma^{2m-1}) - \nu_{-1}(\gamma^{2m-1}) = 2i_1(\gamma) - 1$$

Case 5. $M = -I_2$, the 2×2 identity matrix.

In this case we have

$$i_{-1}(\gamma^{2m+1}) - i_{-1}(\gamma^{2m-1}) - \nu_{-1}(\gamma^{2m-1}) = 2i_1(\gamma) - 2.$$

Case 6. M is conjugate to a matrix $\begin{pmatrix} -1 & b \\ 0 & -1 \end{pmatrix}$ for some b > 0.

In this case we have

$$i_{-1}(\gamma^{2m+1}) - i_{-1}(\gamma^{2m-1}) - \nu_{-1}(\gamma^{2m-1}) = 2i_1(\gamma) - 1$$

Case 7. $M = \begin{pmatrix} \cos \theta & -\sin \theta \\ \sin \theta & \cos \theta \end{pmatrix}$ with some $\theta \in (0, \pi) \cup (\pi, 2\pi)$.

In this case we have

$$\begin{split} &i_{-1}(\gamma^{2m+1})-i_{-1}(\gamma^{2m-1})-\nu_{-1}(\gamma^{2m-1})\\ =&\quad 2(i_1(\gamma)-1)+2E\left(\frac{(2m+1)\theta}{2\pi}+\frac{1}{2}\right)-2E\left(\frac{(2m-1)\theta}{2\pi}+\frac{1}{2}\right)\\ &\quad -\left(2-2\phi\left(\frac{(2m-1)\theta}{2\pi}+\frac{1}{2}\right)\right)\\ \geq&\quad 2(i_1(\gamma)-1). \end{split}$$

Case 8.
$$M = \begin{pmatrix} R(\theta) & b \\ 0 & R(\theta) \end{pmatrix}$$
 with some $\theta \in (0,\pi) \cup (\pi,2\pi)$ and $b = \begin{pmatrix} b_1 & b_2 \\ b_3 & b_4 \end{pmatrix} \in \mathbf{R}^{2\times 2}$, such that $(b_2 - b_3) \sin \theta < 0$.

In this case we have

$$\begin{split} &i_{-1}(\gamma^{2m+1})-i_{-1}(\gamma^{2m-1})-\nu_{-1}(\gamma^{2m-1})\\ =&\quad 2i_{1}(\gamma)+2\phi\left(\frac{(2m+1)\theta}{2\pi}+\frac{1}{2}\right)-2\phi\left(\frac{(2m-1)\theta}{2\pi}+\frac{1}{2}\right)\\ &-\left(2-2\phi\left(\frac{(2m-1)\theta}{2\pi}+\frac{1}{2}\right)\right)\\ \geq&\quad 2i_{1}(\gamma)-2. \end{split}$$

Case 9. $M = \begin{pmatrix} R(\theta) & b \\ 0 & R(\theta) \end{pmatrix}$ with some $\theta \in (0,\pi) \cup (\pi,2\pi)$ and $b = \begin{pmatrix} b_1 & b_2 \\ b_3 & b_4 \end{pmatrix} \in \mathbf{R}^{2\times 2}$, such that $(b_2 - b_3) \sin \theta > 0$.

In this case we have

$$i_{-1}(\gamma^{2m+1}) - i_{-1}(\gamma^{2m-1}) - \nu_{-1}(\gamma^{2m-1})$$

$$= 2i_{1}(\gamma) - \left(2 - 2\phi\left(\frac{(2m-1)\theta}{2\pi} + \frac{1}{2}\right)\right)$$

$$\geq 2i_{1}(\gamma) - 2.$$

Case 10. M is hyperbolic, i.e., $\sigma(M) \cap \mathbf{U} = \emptyset$.

In this case we have

$$i_{-1}(\gamma^{2m+1}) - i_{-1}(\gamma^{2m-1}) - \nu_{-1}(\gamma^{2m-1}) = 2i_1(\gamma).$$

Combining the above cases, we obtain the proposition.

4 Proof of the main theorem

In this section, we give the proof of the main theorem. first we have the following.

Lemma 4.1. Suppose u^{2k-1} is a nonzero critical point of Ψ_a such that u corresponds to $(\tau, y) \in \mathcal{T}_s(\Sigma)$. Then we can find $m \in \mathbf{N}$ such that

$$i(u^{2m+1}) - i(u^{2m-1}) \ge 4. (4.1)$$

Proof. Let $(\tau, y) \in \mathcal{T}_s(\Sigma)$. The fundamental solution $\gamma_y : [0, \tau/2] \to \operatorname{Sp}(2n)$ with $\gamma_y(0) = I_{2n}$ of the linearized Hamiltonian system

$$\dot{w}(t) = JH''(y(t))w(t), \qquad \forall t \in \mathbf{R}, \tag{4.2}$$

is called the associate symplectic path of (τ, y) . Then as in §1.7 of [Eke3], we have

$$\gamma_y(\tau/2) = \begin{pmatrix} -I_2 & 0\\ 0 & C \end{pmatrix} \tag{4.3}$$

in an appropriate coordinate. Then by Lemma 3.1 and Theorem 3.5, we have

$$i(u^{2k-1}) = i_{-1}(\gamma^{2k-1}), \quad \nu(u^{2k-1}) = \nu_{-1}(\gamma^{2k-1}),$$
 (4.4)

for any $k \in \mathbb{N}$. By Theorem 3.10, the matrix $\gamma_y(\tau/2)$ can be connected in $\Omega^0(\gamma_y(\tau/2))$ to a basic form decomposition $M = (-I_2) \diamond M_1 \diamond \cdots \diamond M_l$. Since $n \geq 2$, we may write $M = (-I_2) \diamond M_1 \diamond M'$, where $M' = M_2 \diamond \cdots \diamond M_l$. By the symplectic additivity of indices, cf. [Lon2]-[Lon4], we have

$$i_{-1}(\gamma^{2k-1}) = i_{-1}(\gamma_1^{2k-1}) + i_{-1}(\gamma_2^{2k-1})$$

$$\tag{4.5}$$

where γ_1 and γ_2 are appropriate symplectic paths such that $\gamma_1(\tau/2) = (-I_2) \diamond M_1$ and $\gamma_2(\tau/2) = M'$.

Note that by Theorem 3.5, we have $i_1(\gamma) \geq n$. Now we consider each case as in §3.

Case 1.
$$M_1 = \begin{pmatrix} 1 & b \\ 0 & 1 \end{pmatrix}$$
 for some $b > 0$ or $M_1 = I_2$.

In this case we have

$$i(u^{2m+1}) - i(u^{2m-1}) = i_{-1}(\gamma^{2m+1}) - i_{-1}(\gamma^{2m-1})$$

$$= i_{-1}(\gamma_1^{2m+1}) - i_{-1}(\gamma_1^{2m-1}) + i_{-1}(\gamma_2^{2m+1}) - i_{-1}(\gamma_2^{2m-1})$$

$$\geq 2i_1(\gamma_1) + 2 + 2i_1(\gamma_2) - (2n-4) + \nu_{-1}(\gamma_2^{2m-1})$$

$$\geq 2i_1(\gamma) + 6 - 2n \geq 6.$$

Note that in the above computations, we use (3.29), (3.30), (3.33), Proposition 3.11 and $i_1(\gamma) \geq n$.

Case 2. M is conjugate to a matrix $\begin{pmatrix} 1 & b \\ 0 & 1 \end{pmatrix}$ for some b < 0.

In this case, by (3.31) we have

$$i(u^{2m+1}) - i(u^{2m-1}) = i_{-1}(\gamma^{2m+1}) - i_{-1}(\gamma^{2m-1})$$

$$= i_{-1}(\gamma_1^{2m+1}) - i_{-1}(\gamma_1^{2m-1}) + i_{-1}(\gamma_2^{2m+1}) - i_{-1}(\gamma_2^{2m-1})$$

$$\geq 2i_1(\gamma_1) + 2i_1(\gamma_2) - (2n-4) + \nu_{-1}(\gamma_2^{2m-1})$$

$$\geq 2i_1(\gamma) + 4 - 2n \geq 4.$$

Case 3.
$$M = \begin{pmatrix} -1 & b \\ 0 & -1 \end{pmatrix}$$
 for $b \in \mathbf{R}$.

In this case, by (3.32)-(3.34) we have

$$i(u^{2m+1}) - i(u^{2m-1}) = i_{-1}(\gamma^{2m+1}) - i_{-1}(\gamma^{2m-1})$$

$$= i_{-1}(\gamma_1^{2m+1}) - i_{-1}(\gamma_1^{2m-1}) + i_{-1}(\gamma_2^{2m+1}) - i_{-1}(\gamma_2^{2m-1})$$

$$\geq 2i_1(\gamma_1) + 2i_1(\gamma_2) - (2n-4) + \nu_{-1}(\gamma_2^{2m-1})$$

$$\geq 2i_1(\gamma) + 4 - 2n \geq 4.$$

Case 4. $M = \begin{pmatrix} \cos \theta & -\sin \theta \\ \sin \theta & \cos \theta \end{pmatrix}$ with some $\theta \in (0, \pi) \cup (\pi, 2\pi)$. In this case, by (3.35) we h

$$i(u^{2m+1}) - i(u^{2m-1}) = i_{-1}(\gamma^{2m+1}) - i_{-1}(\gamma^{2m-1})$$

$$= i_{-1}(\gamma_1^{2m+1}) - i_{-1}(\gamma_1^{2m-1}) + i_{-1}(\gamma_2^{2m+1}) - i_{-1}(\gamma_2^{2m-1})$$

$$\geq 2i_1(\gamma_1) - 2 + 2E\left(\frac{(2m+1)\theta}{2\pi} + \frac{1}{2}\right) - 2E\left(\frac{(2m-1)\theta}{2\pi} + \frac{1}{2}\right)$$

$$+2i_1(\gamma_2) - (2n-4) + \nu_{-1}(\gamma_2^{2m-1})$$

$$\geq 2i_1(\gamma) + 4 - 2n \geq 4$$

provided we choose m such that $E\left(\frac{(2m+1)\theta}{2\pi}+\frac{1}{2}\right)-E\left(\frac{(2m-1)\theta}{2\pi}+\frac{1}{2}\right)\geq 1$.

Case 5. $M=\begin{pmatrix}R(\theta)&b\\0&R(\theta)\end{pmatrix}$ with some $\theta\in(0,\pi)\cup(\pi,2\pi)$ and $b=\begin{pmatrix}b_1&b_2\\b_3&b_4\end{pmatrix}\in\mathbf{R}^{2\times2}$, such that $(b_2 - b_3) \sin \theta < 0$.

In this case, by (3.36) we have

$$\begin{split} &i(u^{2m+1})-i(u^{2m-1})=i_{-1}(\gamma^{2m+1})-i_{-1}(\gamma^{2m-1})\\ =&\quad i_{-1}(\gamma_1^{2m+1})-i_{-1}(\gamma_1^{2m-1})+i_{-1}(\gamma_2^{2m+1})-i_{-1}(\gamma_2^{2m-1})\\ \geq&\quad 2i_1(\gamma_1)+2\varphi\left(\frac{(2m+1)\theta}{2\pi}+\frac{1}{2}\right)-2\varphi\left(\frac{(2m-1)\theta}{2\pi}+\frac{1}{2}\right)\\ &\quad +2i_1(\gamma_2)-(2n-6)+\nu_{-1}(\gamma_2^{2m-1})\\ \geq&\quad 2i_1(\gamma)+4-2n\geq 4. \end{split}$$

Case 6. $M = \begin{pmatrix} R(\theta) & b \\ 0 & R(\theta) \end{pmatrix}$ with some $\theta \in (0, \pi) \cup (\pi, 2\pi)$ and $b = \begin{pmatrix} b_1 & b_2 \\ b_3 & b_4 \end{pmatrix} \in \mathbf{R}^{2 \times 2}$, such that $(b_2-b_3)\sin\theta>0$

In this case, by (3.37) we have

$$i(u^{2m+1}) - i(u^{2m-1}) = i_{-1}(\gamma^{2m+1}) - i_{-1}(\gamma^{2m-1})$$

$$= i_{-1}(\gamma_1^{2m+1}) - i_{-1}(\gamma_1^{2m-1}) + i_{-1}(\gamma_2^{2m+1}) - i_{-1}(\gamma_2^{2m-1})$$

$$\geq 2i_1(\gamma_1) + 2i_1(\gamma_2) - (2n-6) + \nu_{-1}(\gamma_2^{2m-1})$$

$$\geq 2i_1(\gamma) + 6 - 2n \geq 6.$$

Case 7. M is hyperbolic, i.e., $\sigma(M) \cap \mathbf{U} = \emptyset$.

In this case, by (3.38) we have

$$i(u^{2m+1}) - i(u^{2m-1}) = i_{-1}(\gamma^{2m+1}) - i_{-1}(\gamma^{2m-1})$$

$$= i_{-1}(\gamma_1^{2m+1}) - i_{-1}(\gamma_1^{2m-1}) + i_{-1}(\gamma_2^{2m+1}) - i_{-1}(\gamma_2^{2m-1})$$

$$\geq 2i_1(\gamma_1) + 2i_1(\gamma_2) - (2n-4) + \nu_{-1}(\gamma_2^{2m-1})$$

$$\geq 2i_1(\gamma) + 4 - 2n \geq 4.$$

Combining all the above cases, we obtain the lemma.

Proof of Theorem 1.1. We prove by contraction. Assume $\mathcal{T}_s(\Sigma) = \{(\tau, y)\}$. Suppose u^{2m-1} is a nonzero critical point of Ψ_a such that u corresponds to $(\tau, y) \in \mathcal{T}_s(\Sigma)$. By Lemma 4.1, we may assume $i(u^{2m+1}) - i(u^{2m-1}) \geq 4$. The index interval of (τ, y) at 2m-1 is defined to be $\mathcal{G}_{2m-1} = (i(u^{2m-3}) + \nu(u^{2m-3}) - 1, i(u^{2m+1}))$. Note that by Proposition 3.11 and $i_1(\gamma) \geq n$, we have $i(u^{2m-3}) + \nu(u^{2m-3}) \leq i(u^{2m-1})$. Thus we have $(i(u^{2m-1}) - 1, i(u^{2m+1})) \subset \mathcal{G}_{2m-1}$. Hence we can find two distinct even integers $2T_1, 2T_2 \in \mathcal{G}_{2m-1}$. Let c_{T_1+1} and c_{T_2+1} be the two critical values of Ψ_a found by Proposition 2.16. Then we have $c_{T_1+1} \neq c_{T_2+1}$ since ${}^{\#}\mathcal{T}_s(\Sigma) < \infty$. By Proposition 2.17, we have

$$\Psi_a(u^{2m_1-1}) = c_{T_1+1}, \qquad i(u^{2m_1}) \le 2T_1 \le i(u^{2m_1-1}) + \nu(u^{2m_1-1}) - 1,
\Psi_a(u^{2m_2-1}) = c_{T_2+1}, \qquad i(u^{2m_2}) \le 2T_2 \le i(u^{2m_2-1}) + \nu(u^{2m_2-1}) - 1, \tag{4.6}$$

for some $m_1, m_2 \in \mathbf{N}$. On the other hand, we must have $m_1 = m_2$ by Proposition 3.11. Thus we have $c_{T_1+1} = c_{T_2+1}$. This contradiction proves the theorem.

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